



The Case of Neumann, Robin, and Periodic Lateral Conditions for the Semi-infinite Generalized Graetz Problem and Applications

Valentin Debarnot, Jérôme Fehrenbach, Frédéric de Gournay, Léo Martire

► To cite this version:

Valentin Debarnot, Jérôme Fehrenbach, Frédéric de Gournay, Léo Martire. The Case of Neumann, Robin, and Periodic Lateral Conditions for the Semi-infinite Generalized Graetz Problem and Applications. SIAM Journal on Applied Mathematics, 2018, 78 (4), pp.2227 - 2251. 10.1137/17M1157507 . hal-01880777

HAL Id: hal-01880777

<https://hal.science/hal-01880777>

Submitted on 25 Sep 2018

HAL is a multi-disciplinary open access archive for the deposit and dissemination of scientific research documents, whether they are published or not. The documents may come from teaching and research institutions in France or abroad, or from public or private research centers.

L'archive ouverte pluridisciplinaire **HAL**, est destinée au dépôt et à la diffusion de documents scientifiques de niveau recherche, publiés ou non, émanant des établissements d'enseignement et de recherche français ou étrangers, des laboratoires publics ou privés.

THE CASE OF NEUMANN, ROBIN AND PERIODIC LATERAL CONDITION FOR THE SEMI INFINITE GENERALIZED GRAETZ PROBLEM AND APPLICATIONS.

VALENTION DEBARNOT, JÉRÔME FEHRENBACH, FRÉDÉRIC DE GOURNAY, LÉO
MARTIRE *

Abstract. The Graetz problem is a convection-diffusion equation in a pipe invariant along a direction. The contribution of the present work is to propose a mathematical analysis of the Neumann, Robin and periodic boundary condition on the boundary of a semi-infinite pipe. The solution in the 3D space of the original problem is reduced to eigenproblems in the 2D section of the pipe. The set of solutions is described, its structure depends on the type of boundary condition and of the sign of the total flow of the fluid. This analysis is the cornerstone of numerical methods to solve Graetz problem in finite pipes, semi infinite pipes and exchangers of arbitrary cross section. Numerical test-cases illustrate the capabilities of these methods to provide solutions in various configurations.

1. Introduction.

1.1. Context. The seminal work of Graetz in the late 19th century addressed a stationary convection-diffusion problem inside an axi-symmetrical cylindrical pipe [5], where the regime was supposed to be convection-dominated which means that the longitudinal diffusion was neglected. It was the first contribution to the modelling of convective transport coupled with diffusion, with important applications nowadays as the parallel convective exchangers involved in heating or cooling systems [16], haemodialysis [1], and heat exchangers [7]. The first extension to the Graetz problem, known as the “extended Graetz problem” takes into account longitudinal diffusion [10, 3, 18, 9]. Papoutsakis *et al.* in [12, 11] introduced a symmetric operator acting on a two-components space that solves the extended Graetz problem in axi-symmetrical configurations. The so-called “conjugated Graetz problem” where multiple solid or fluid phases are taken into account was proposed in [13, 14] in the case of an axi-symmetrical configuration. These successive models aimed at taking into account more and more complex and realistic situations, and when only axi-symmetrical configurations were considered the equations boiled down to one-dimensional problems. The adaptation to parallel plates heat exchangers of these one-dimensional models, together with a parametric study was proposed in [6]. The reader may also consult [2] for a review on the conjugated Graetz problem.

The work on non-axisymmetrical configurations was initiated in [15] where the operator was proved to be self-adjoint with compact resolvent when Dirichlet boundary conditions are applied on the boundary of the domain. In the case of a single fluid stream the negative eigenvalues correspond to downstream propagation, and positive eigenvalues to upstream propagation. The main novelty was that arbitrary geometries were addressed, and a detailed mathematical analysis of the Dirichlet problem was proposed. The authors of [15] coined this problem as the “generalized Graetz problem”. Numerical methods for the approximation of this operator and error estimates were provided in [4].

The objective of the present work is to extend the work of [4] and provide explicit methods with general lateral boundary conditions, beyond the Dirichlet case. The cross section of the domain has an arbitrary geometry and can incorporate different

* Institut de Mathématiques de Toulouse (UMR 5219), Université de Toulouse, CNRS UPS, F-31062 Toulouse, France F-31077 Toulouse, France INSA, F-31077 Toulouse, France (frederic@degournay.fr), (jerome.fehrenbach@math.univ-toulouse.fr)

fluid domains, possibly with opposite signs of the velocity. The lateral boundary conditions that we address can be Dirichlet, Neumann, Robin, periodic or a mixture of these different cases on different parts of the boundary. The periodic boundary conditions with rectangular or hexagonal cell are adapted to the analysis of micro-exchangers, where a design pattern is repeated.

1.2. Setting. In convection-dominated heat or mass transfer, we address the generalized Graetz problem which occurs in a cylinder of arbitrary section Ω and of length I , possibly $I = \mathbb{R}^+$, see Figure 1.1. The diffusion coefficient is supposed to be invariant by translation along e_z the axis of the cylinder. Similarly, the velocity vector v is supposed to be oriented in the direction of the axis of the cylinder, that is $v = he_z$ with $h \in L^\infty(\Omega)$. The equation for the temperature T inside the domain is then

$$c\partial_{zz}T + \operatorname{div}(\sigma\nabla T) - h\partial_z T = 0 \text{ on } \Omega \times I, \quad (\text{E})$$

with diffusion coefficients $c, \sigma > 0$ bounded in Ω with bounded inverse. In (E), it is implicitly supposed that the heat capacity C and the density ρ of the fluid satisfy $\rho C = 1$. If one has to handle several fluids with different physical properties, the choice of an adequate normalization leads to (E). In this case, at each point $x \in \Omega$, $h(x)$ represents the velocity multiplied by $\rho(x)C(x)$. For simplicity reasons, we refer in the sequel to h as the *velocity*.

The lateral homogeneous boundary conditions (LBC) may be of Neumann, Dirichlet, Robin and periodic type, respectively on $\Gamma_N, \Gamma_D, \Gamma_R, \Gamma_\# \subset \partial\Omega$ given by

$$\begin{cases} \sigma\nabla T \cdot n = 0 \text{ on } \Gamma_N \times I: \text{Neumann, and/or} \\ T = 0 \text{ on } \Gamma_D \times I: \text{Dirichlet, and/or} \\ \sigma\nabla T \cdot n + aT = 0 \text{ on } \Gamma_R \times I: \text{Robin, and/or} \\ T \text{ is periodic on } \Gamma_\# \times I: \text{periodic} \end{cases} \quad (\text{LBC})$$

where $a > 0$ in the Robin condition, and $\Gamma_\#$ must be tailored to support periodic conditions (e.g. Ω is the unit square, $\Gamma_\# = (\{x = 0\} \cup \{x = 1\}) \cap \partial\Omega$ and the boundary condition is $T(0, y) = T(1, y)$). As usual, the Γ 's involved in the definition of the boundary condition must form a partition of $\partial\Omega$. Note that the Neumann (resp. Dirichlet) boundary conditions are degenerate cases of the Robin condition corresponding to $a = 0$ (resp. $a = +\infty$). The Inlet/Outlet boundary condition (I/OBC) is of Dirichlet and/or of Neumann type and is given by

$$T = T_D \text{ on } \Omega_D \text{ and } \partial_z T = S_N \text{ on } \Omega_N \text{ with } \Omega_D \cup \Omega_N = \Omega \times \partial I. \quad (\text{I/OBC})$$

In the case $I = \mathbb{R}^+$, we intentionally stay vague about the definition of ∂I , it is one of the results of this work to determine whether an (I/OBC) is needed on $z = +\infty$.

A more realistic model in regimes of high velocities takes into account a viscosity term, see e.g. [8] where a study in a microchannel including viscous effects and longitudinal conduction is performed. Our approach can also account for viscosity, the details are presented in Section 2.3.

1.3. Lax-Milgram. Note that the equation (E) is an elliptic equation with an additional convective term. It is possible to use Lax-Milgram's theorem [17] under the hypothesis that the Inlet/Outlet boundary condition is Dirichlet in the region where the flow is incoming. More precisely:

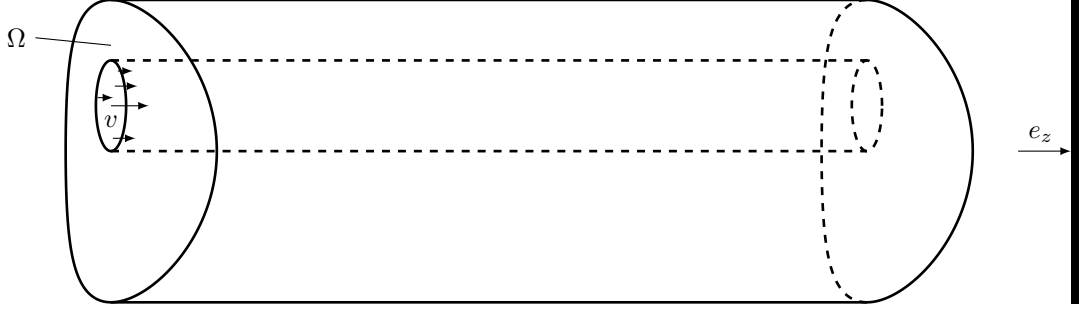


FIGURE 1.1. The domain $\Omega \times I$ where the Graetz problem is posed.

PROPOSITION 1. Let $I = [z_1, z_2]$ and $\omega_{\pm} = \{x \text{ s.t. } \pm h(x) > 0\}$. If

$$\omega_+ \times \{z_1\} \subset \Omega_D \text{ and } \omega_- \times \{z_2\} \subset \Omega_D$$

and if T_D and S_N are regular enough, then there exists a unique solution to (E) with the boundary conditions (LBC) and (I/OBC).

The proof is only sketched here for the sake of completeness. Denote \mathcal{X} the natural space of elements where the solution is sought, that is

$$\mathcal{X} = \{T \in H^1(\Omega \times I) \text{ s.t. } T = 0 \text{ on } (\Gamma_D \times I) \cup \Omega_D \text{ and } T \text{ periodic on } \Gamma_{\#} \times I\}.$$

Non-homogeneous Dirichlet boundary conditions of (I/OBC) are solved using a lift of T_D , still denoted T_D that satisfies the lateral boundary conditions (LBC) with $\partial_z T_D = 0$ on Ω_N and denote

$$f_D = c \partial_{zz} T_D + \operatorname{div}(\sigma \nabla T_D) - h \partial_z T_D.$$

The change of unknown $\tilde{T} = T - T_D$ where T solves (E) and (LBC), leads to the following variational formulation: find $\tilde{T} \in \mathcal{X}$ such that for every $\phi \in \mathcal{X}$:

$$\underbrace{\int_{\Omega \times I} c \partial_z \tilde{T} \partial_z \phi + \sigma \nabla \tilde{T} \cdot \nabla \phi + h \partial_z \tilde{T} \phi + \int_{\Gamma_R} a \tilde{T} \phi}_{b(\tilde{T}, \phi)} + \underbrace{\int_{\Omega_N} S_N \phi - \int_{\Omega \times I} f_D \phi}_{\ell(\phi)} = 0.$$

The term $b(T, \phi)$ is bilinear in (T, ϕ) and continuous for the standard norm of \mathcal{X} , the term $\ell(\phi)$ is linear continuous if T_D and S_N are regular enough. It remains to study the coercivity of b .

$$b(T, T) = \int_{\Omega \times I} c \partial_z T \cdot \partial_z T + \sigma \nabla T \cdot \nabla T + h \partial_z T \cdot T = \int_{\Omega \times I} \|\nabla_{3D} T\|_{\kappa}^2 + \frac{1}{2} \int_{\Omega \times I} h \partial_z (T^2),$$

where κ is a positive matrix with diagonal entries (σ, σ, c) in the basis (e_x, e_y, e_z) . The first term is coercive. The second term is

$$\frac{1}{2} \int_{\Omega} h T^2|_{z=z_2} - \frac{1}{2} \int_{\Omega \times \{z_2\}} h T^2 - \frac{1}{2} \int_{\Omega \times \{z_1\}} h T^2.$$

It is nonnegative for all $T \in \mathcal{X}$ if and only if the Inlet/Outlet condition is of Dirichlet type at the boundary where the flow is entering the domain ($z = z_1$ if $h > 0$, and $z = z_2$ if $h < 0$).

1.4. Presentation of the paper. The objective of the present paper is to provide a general framework that allows to solve (E) with any type of boundary condition beyond the case where Lax-Milgram's theorem can be used. Section 2 details the notation and the main properties of the operator involved in the solution, as well as the modifications required to take into account a viscosity term. The main results, namely Theorems 4 and 5, are detailed in Section 3, their proof is postponed to the Appendix. In Section 4 we solve the problem in a semi-infinite domain and show that depending on the case the temperature at infinity T_∞ can either be a free parameter of the problem or be imposed by the other condition. In Section 5 we address the case of a domain of finite length, and numerical strategies are detailed in the different cases depending on the lateral boundary condition and on the Inlet/Outlet condition. Test cases are presented in Section 6.

2. State of the art and position of the problem. The equation (E) may be interpreted as an evolution equation in the variable z if it is cast into

$$\partial_z \begin{pmatrix} \partial_z T \\ T \end{pmatrix} = \mathcal{A} \begin{pmatrix} \partial_z T \\ T \end{pmatrix} \text{ on } \Omega \times I, \text{ with } \mathcal{A} \begin{pmatrix} u \\ s \end{pmatrix} = \begin{pmatrix} hc^{-1}u - c^{-1} \operatorname{div} \sigma \nabla s \\ u \end{pmatrix}. \quad (2.1)$$

The goal of this section is to guide the reader to the analysis of (2.1) that was proposed in [4], to enlarge the frame to Neumann and periodic lateral boundary condition, and to define the notation and state the results that will be used in the sequel. Since \mathcal{A} is a symmetric operator with a compact resolvent, classical eigendecomposition leads to an explicit representation of the solution of (2.1) in the basis of eigenvectors (see e.g. [17]).

DEFINITION 2. *We say that “the constants are not controlled” when $\Gamma_D \cup \Gamma_R = \emptyset$, in other words when there is no Dirichlet or Robin condition on the lateral part of the boundary of the domain. The case where the constants are not controlled and in addition $\int_\Omega h = 0$ is called the “balanced case”.*

From the engineering point of view, the balanced case is a special instance of counter-current configuration, where the integral of the velocities in the 2 directions have the same magnitude and the boundary of the domain Ω is perfectly insulating or periodic. As we prove in this section, the case where the constants are not controlled is a case where the constants are a solution of (E) and the balanced case is a case where \mathcal{A} admits a non-trivial kernel.

2.1. Study of the operator \mathcal{A} . In this section we detail the Hilbert space, the scalar product, the kernel, range and pseudo-inverse of the symmetric operator \mathcal{A} .

Hilbert space and scalar product. First, introduce the space H that encodes the lateral boundary condition. When the constants are controlled define:

$$H = \{s \in H^1(\Omega), \text{ such that } s = 0 \text{ on } \Gamma_D \text{ and } s \text{ periodic on } \Gamma_\# \}.$$

If there is no Dirichlet or Robin boundary condition, hence no control on the constants, quotient by the constants and define:

$$H = \{s \in H^1(\Omega)/\mathbb{R}, \text{ such that } s \text{ periodic on } \Gamma_\# \}.$$

Then, define the Hilbert space \mathcal{H} as

$$\mathcal{H} = \{(u, s) \mid u \in L^2(\Omega), s \in H\}$$

which is endowed with the scalar product:

$$((u, s) | (u', s'))_{\mathcal{H}} = \int_{\Omega} cuu' + \sigma \nabla s \cdot \nabla s' + \int_{\Gamma_R} ass'.$$

101 The crucial step in showing that \mathcal{H} is a Hilbert space is to show that the scalar product
 102 is definite. Setting $((u, s) | (u, s))_{\mathcal{H}} = 0$ immediatly gives $u = 0$ and $\nabla s = 0$, hence s
 103 is a constant. If the constants are controlled, then $\Gamma_D \cup \Gamma_R \neq \emptyset$ and $s = 0$, whereas
 104 if the constants are not controlled then s is a constant and $s = 0$ in H .

The domain of the operator \mathcal{A} is:

$$\mathcal{D}(\mathcal{A}) = \{(u, s) \in \mathcal{H}, u \in H^1(\Omega), \operatorname{div}(\sigma \nabla s) \in L^2(\Omega) + \text{boundary conditions (LBC)}\},$$

105 where the boundary conditions are $u \in H$, and $\sigma \nabla s \cdot n$ is equal to 0 on Γ_N , is equal
 106 to $-as$ on Γ_R and is periodic on Γ_{\sharp} . On $\mathcal{D}(\mathcal{A})$, the operator is symmetric as we prove
 107 now. Let $\phi = (u, s)$ and $\phi' = (u', s') \in \mathcal{D}(\mathcal{A})$:

$$\begin{aligned} (\mathcal{A}\phi | \phi')_{\mathcal{H}} &= \int_{\Omega} (hu - \operatorname{div} \sigma \nabla s)u' + \sigma \nabla u \cdot \nabla s' + \int_{\Gamma_R} aus' \\ &= \int_{\Omega} hu u' + \sigma \nabla u \cdot \nabla s' + \underbrace{\sigma \nabla u' \cdot \nabla s + \int_{\partial \Omega} (-\sigma \nabla s \cdot n)u' + \int_{\Gamma_R} aus'}_{(1)}, \end{aligned}$$

108 and the term (1) is symmetric thanks to (LBC) on $\mathcal{D}(\mathcal{A})$.

Inverse of the Laplacian Define the inverse of the Laplace operator as:

$$u = \Delta_{\sigma}^{-1} f \text{ iff } \begin{cases} \operatorname{div}(\sigma \nabla u) = f, \text{ and} \\ u \in H \\ + \text{boundary conditions,} \end{cases}$$

109 where the boundary conditions are $\sigma \nabla u \cdot n = 0$ on Γ_N and $\sigma \nabla u \cdot n + au = 0$ on Γ_R .
 110 If the constants are controlled, then Δ_{σ}^{-1} is well defined on $L^2(\Omega)$, whereas if there
 111 is only Neumann or periodic boundary conditions (no control of the constants), the
 112 operator Δ_{σ}^{-1} is only defined if $f \in L_m^2(\Omega)$, the subspace of $L^2(\Omega)$ with null average.

Kernel of \mathcal{A} Following from the definition of \mathcal{A} in (2.1), the kernel of \mathcal{A} is the set of (u, s) in $\mathcal{D}(\mathcal{A})$ such that

$$u = 0 \text{ in } H \text{ and } hu - \operatorname{div}(\sigma \nabla s) = 0.$$

When the constants are controlled, both u and s are then equal to 0. When the constants are not controlled, since u is a constant, then $s = u \Delta_{\sigma}^{-1} h$ in Ω which admits a solution if and only if $\int_{\Omega} h = 0$. To summarize the kernel of \mathcal{A} is:

$$\mathcal{K}(\mathcal{A}) = \begin{cases} \operatorname{Vect}(\phi_0 = (1, \Delta_{\sigma}^{-1} h)) \text{ in the balanced case ,} \\ \{0\} \text{ in the other cases} \end{cases}$$

Range and inverse of \mathcal{A} The range of \mathcal{A} , denoted $\mathcal{R}(\mathcal{A})$, is defined as the orthogonal of $\mathcal{K}(\mathcal{A})$ in \mathcal{H} and the inverse of \mathcal{A} is an operator from $\mathcal{R}(\mathcal{A})$ to $\mathcal{D}(\mathcal{A})$, defined as follows:

$$\forall \phi = (u, s) \in \mathcal{R}(\mathcal{A}), \mathcal{A}^{-1} \phi = \begin{cases} (s, \Delta_{\sigma}^{-1}(hs - cu)) \text{ if the constants are controlled} \\ (s + k, \Delta_{\sigma}^{-1}(hs - cu + hk)), k \in \mathbb{R} \text{ if not.} \end{cases}.$$

When the constants are not controlled, the constant $k \in \mathbb{R}$ is chosen so that

$$\begin{cases} \int_{\Omega} hs - cu + hk = 0 & \text{in the non-balanced case,} \\ (\mathcal{A}^{-1}\phi, \phi_0)_{\mathcal{H}} = 0 & \text{in the balanced case.} \end{cases}$$

It is easily checked that for all $\phi \in \mathcal{R}(\mathcal{A})$, $\mathcal{A}^{-1}\phi \in \mathcal{D}(\mathcal{A})$ and that $\mathcal{A}\mathcal{A}^{-1}\phi = \phi$. The operator \mathcal{A}^{-1} is then symmetric (as a consequence of the symmetry of \mathcal{A}). Note also that in the balanced case, one can also write:

$$\forall \phi = (u, s) \in \mathcal{R}(\mathcal{A}), \quad \mathcal{A}^{-1}\phi = (s, \Delta_{\sigma}^{-1}(hs - cu)) + k\phi_0.$$

Eigenvalue decomposition of \mathcal{A} The operator \mathcal{A}^{-1} is a compact self-adjoint operator on $\mathcal{R}(\mathcal{A})$. To prove this let $\phi_n = (u_n, s_n)$ be a bounded sequence in \mathcal{H} . Then up to a subsequence it is a weakly convergent sequence and s_n converges strongly in $L^2(\Omega)$. Using the fact that Δ_{σ}^{-1} is a compact operator from L^2 to H finishes the proof. We denote by λ_i the non-zero ordered eigenvalues of \mathcal{A} and by $\phi_i = (U_i, \lambda_i^{-1}U_i)$ the corresponding eigenvectors. By convention, λ_i is of the sign of i so that

$$-\infty \leftarrow \lambda_{-n} \leq \lambda_{-n-1} \leq \dots \leq \lambda_{-1} < 0 < \lambda_1 \leq \dots \leq \lambda_{n-1} \leq \lambda_n \rightarrow +\infty.$$

113 In the balanced case, we add to the family $(\phi_i)_i$ the vector $\phi_0 = (1, \Delta_{\sigma}^{-1}h)$, so
114 that the Hilbert space \mathcal{H} is the space spanned by the eigenvectors $(\phi_i)_{i \in \mathbb{Z}}$.

2.2. Solution of the evolution equation. The diagonalization of the operator \mathcal{A} allows to solve the evolution equation (E):

$$c\partial_{zz}T + \operatorname{div}(\sigma \nabla T) - h\partial_z T = 0 \text{ on } \Omega \times I. \quad (\text{E})$$

Let $T \in C^1(I, L^2(\Omega)) \cap C^0(I, H)$ be a solution of this equation with corresponding lateral boundary conditions (LBC). If we denote $\phi : z \mapsto (\partial_z T(z), T(z))$ in $C^0(\mathcal{H})$, then the equation (E) is equivalent to $\partial_z \phi = \mathcal{A}\phi$, and the solution ϕ is given by

$$\phi(z) = \sum_{i \in \mathbb{Z}} \frac{(\phi(0)|\phi_i)_{\mathcal{H}}}{\|\phi_i\|_{\mathcal{H}}^2} e^{\lambda_i z} \phi_i. \quad (2.2)$$

One can either identify the first coordinate and integrate w.r.t. z or identify the second coordinate and denote

$$\psi = \sum_{i \in \mathbb{Z}^*} \lambda_i^{-1} \frac{(\phi(0)|\phi_i)_{\mathcal{H}}}{\|\phi_i\|_{\mathcal{H}}^2} \phi_i,$$

115 to obtain

$$T(z) = \sum_{i \in \mathbb{Z}^*} (\psi|\phi_i)_{\mathcal{H}} U_i e^{\lambda_i z} \text{ if the constants are controlled, i.e. } \Gamma_D \cup \Gamma_R \neq \emptyset,$$

$$T(z) = \sum_{i \in \mathbb{Z}^*} (\psi|\phi_i)_{\mathcal{H}} U_i e^{\lambda_i z} + a_0 \text{ with } a_0 \in \mathbb{R} \text{ if } \Gamma_D \cup \Gamma_R = \emptyset \text{ and } \int_{\Omega} h \neq 0,$$

$$T(z) = \sum_{i \in \mathbb{Z}^*} (\psi|\phi_i)_{\mathcal{H}} U_i e^{\lambda_i z} + a_0 + a_1(z + \Delta_{\sigma}^{-1}h) \text{ in the balanced case}$$

$$\text{with } a_0 \in \mathbb{R} \text{ and } a_1 = \frac{(\phi(0)|\phi_0)_{\mathcal{H}}}{\|\phi_0\|^2}.$$

If $\partial_z T$ and T are given at $z = 0$ such that $\phi(0) = (\partial_z T(0), T(0))$ belongs to \mathcal{H} , then ψ is uniquely determined. Moreover the constant a_0 is also determined by $T|_{z=0}$ (and also a_1 in the balanced case). We stress that this solution may not be defined everywhere, indeed the series on the right-hand side of (2.2) has to be convergent in some sense and the convergence of the series for $z = 0$ is not sufficient to ensure the convergence for $z \neq 0$ due to the multiplication by $e^{\lambda_i z}$ for non-zero λ_i 's. The set of initial datum ϕ that allows this series to exist is known as the set of compatible initial condition for the Cauchy problem.

2.3. Including a viscous term. Let us consider the following modification of the equation (E) where a viscous term is added:

$$c\partial_{zz}T + \operatorname{div}(\sigma\nabla T) - h\partial_z T = \mu|\nabla h|^2. \quad (2.3)$$

PROPOSITION 3. *Let T be the solution of the Graetz equation with viscosity (2.3). Then there exists an explicit change of unknown function that transforms the problem with viscosity into a problem without viscosity of the form (E). Therefore the solution of the problem (2.3) reduces to the solution of the original problem (E).*

Proof. Once a particular solution \tilde{T} is found, the change of variable $\hat{T} = T - \tilde{T}$ transforms by linearity the problem with viscosity (2.3) into the problem without viscosity. We distinguish different cases, depending on if the constants are controlled or not, and in the case the constants are not controlled we treat separately the non-balanced and the balanced case. In each case we provide an explicit particular solution \tilde{T} .

a) If the constants are controlled, a particular solution is given by

$$\tilde{T} = \Delta_\sigma^{-1}(\mu|\nabla h|^2).$$

b) If the constants are not controlled, in the non-balanced case

$$\tilde{T} = \alpha z + \Delta_\sigma^{-1}(\mu|\nabla h|^2 + \alpha h),$$

where $\alpha \in \mathbb{R}$ satisfies

$$\int_\Omega (\mu|\nabla h|^2 + \alpha h) = 0.$$

c) If the constants are not controlled, in the balanced case, the particular solution is given by $\tilde{T} = \alpha(\frac{z^2}{2} + z\Delta_\sigma^{-1}h) + \Delta_\sigma^{-1}\gamma$, with $\alpha \in \mathbb{R}$ and $\gamma \in L^2(\Omega)$ such that:

$$\begin{cases} \alpha \left(\int_\Omega c - h\Delta_\sigma^{-1}h \right) = \int_\Omega \mu|\nabla h|^2 \\ \gamma = \mu|\nabla h|^2 - \alpha(c - h\Delta_\sigma^{-1}h) \end{cases} \quad (2.4)$$

The choice of α ensures that γ has zero average so that $\Delta_\sigma^{-1}\gamma$ is well defined. Note that α is well defined since

$$\int_\Omega c - h\Delta_\sigma^{-1}h = \int_\Omega c + \int_\Omega \sigma|\nabla h|^2 > 0.$$

Note that the last term is equal to $\|\phi_0\|_{\mathcal{H}}^2$. \square

3. Main decomposition theorem. In this section, the decomposition of a temperature field on the non-positive eigenspace is studied. The result stated in Theorem 5 considers different cases depending on the control of constants and the sign of the total flow.

3.1. Notation and statement of the problem. The $D(\mathcal{A}^\alpha)$ norm or “ α -norm” in short is defined by

$$\|\phi\|_\alpha^2 = \sum_{i \in \mathbb{Z}} \lambda_i^{2\alpha} \frac{(\phi_i|\phi)_{\mathcal{H}}^2}{\|\phi_i\|_{\mathcal{H}}^2} \quad \forall \phi \in \mathcal{H}.$$

The space $D(\mathcal{A}^\alpha)$ is the set of $\phi \in \mathcal{R}(\mathcal{A})$ whose α -norm is $< +\infty$. It is easy to check that $D(\mathcal{A}^1) = D(\mathcal{A})$ and that $D(\mathcal{A}^0) = \mathcal{R}(\mathcal{A})$. Define P an orthogonal projection on \mathcal{H} as:

$$\forall \phi = (u, s) \in \mathcal{H}, \quad P\phi = (u, 0). \quad (3.1)$$

For any I subset of \mathbb{Z} define π_I the orthogonal projection

$$\pi_I \phi = \sum_{i \in I} \frac{(\phi_i|\phi)_{\mathcal{H}}}{\|\phi_i\|_{\mathcal{H}}^2} \phi_i \quad (3.2)$$

We denote $\pi_+ = \pi_{\mathbb{N}^*}$, $\pi_- = \pi_{-\mathbb{N}^*}$, $\pi_0 = \pi_{\{0\}}$ and $\mathcal{R}(\pi_I) = \pi_I(\mathcal{H})$.

The problem of decomposition of a temperature field on the non-positive eigenspace is stated as follows:

For any $\phi \in \mathcal{H}$, find ψ such that

$$P\psi = P\phi \text{ and } \pi_+\psi = 0. \quad (3.3)$$

A similar problem of decomposition on the non-negative eigenspace is obtained by replacing π_+ by π_- . All the results of the present section have a counterpart obtained by changing the sign of z .

3.2. Necessary and sufficient condition. In order to tackle problem (3.3), we first consider the following related problem:

$$\text{Find } \psi \in \mathcal{R}(\pi_-) \text{ such that } \pi_- P \pi_- \psi = \pi_- P \phi. \quad (3.4)$$

Indeed if ψ solves (3.3), then multiplying the equation by π_- and assuming that the kernel of \mathcal{A} is reduced to the nullspace (which is true except in the balanced case), one derives equation (3.4). Such a problem admits a unique solution, given by the following theorem:

THEOREM 4. *The operator $\pi_- P \pi_-$ is invertible on $\mathcal{R}(\pi_-)$. Define B_- the self-adjoint operator of \mathcal{H} as*

$$B_- \phi = \pi_- (\pi_- P \pi_-)^{-1} \pi_- \phi.$$

Moreover it holds

$$\|B_- \phi\|_{\mathcal{H}} \leq C \|\pi_- \phi\|_{\mathcal{H}} \quad \text{and} \quad \|B_- \phi\|_{1/2} \leq C \|\pi_- \phi\|_{1/2}.$$

One can similarly define an operator B_+ , obtained by replacing π_- by π_+ .

The result is proved in [4] for the full-Dirichlet case, that is $\Gamma_D = \partial\Omega$. The proof can be adapted without major changes to the case in consideration. It is reproduced

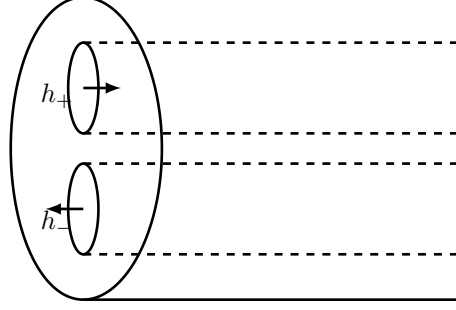


FIGURE 4.1. Example of a semi-infinite cylinder, with two fluid domains.

in Appendix B for the convenience of the reader. Problem (3.3) is then solved in the next theorem.

THEOREM 5. Let $\phi \in \mathcal{H}$, define $\Phi = (1, 0) \in \mathcal{H}$ and consider problem (3.3) of finding ψ a solution of

$$P\psi = P\phi \text{ and } \pi_+\psi = 0. \quad (3.3)$$

- If the constants are not controlled and $\int_{\Omega} h > 0$, there exists a solution if and only if

$$(\Phi - PB_-\Phi|\phi)_{\mathcal{H}} = 0.$$

In this case, the solution is unique and given by $\psi = B_-P\phi$.

- If the constants are not controlled and $\int_{\Omega} h = 0$ (balanced case), then $(\Phi - PB_-\Phi|\Phi)_{\mathcal{H}} \neq 0$ and there exists a unique solution given by

$$\psi = B_-P\phi + \frac{(\Phi - PB_-\Phi|\phi)_{\mathcal{H}}}{(\Phi - PB_-\Phi|\Phi)_{\mathcal{H}}} (B_-\Phi - \phi_0).$$

- In every other case then $\psi = B_-P\phi$ is the unique solution.

The proof of this result is given in Appendix C.

4. Resolution of the semi-infinite problem. In the semi-infinite problem, the equation is set on the cylinder $\Omega \times \mathbb{R}^+$, see Figure 4.1. The equation (E) becomes:

$$\partial_{zz}T + \Delta T + h\partial_z T = 0 \text{ on } \Omega \times \mathbb{R}^+.$$

In this section we address different cases depending on the type of the Inlet/Outlet condition, namely either Dirichlet or Neumann.

In order to ensure uniqueness of the solution, we add the extra hypothesis that the temperature does not grow exponentially. We will say that the temperature has subexponential growth if and only if for every $\lambda > 0$, then $T(z) = o(e^{\lambda z})$ as z goes to $+\infty$.

4.1. Semi-infinite problem, Dirichlet Inlet/Outlet Condition. We consider the Dirichlet (I/OBC) condition:

$$T|_{z=0} = T_0 \quad \text{on } \Omega. \quad (4.1)$$

Denote $\phi_{\mathcal{D}} = (T_0, 0) \in \mathcal{H}$.

PROPOSITION 6. Consider the Graetz problem (E) on the semi-infinite cylinder $\Omega \times [0, +\infty)$, with subexponential growth together with Dirichlet Inlet/Outlet condition (4.1).

a) If the constants are controlled, then there exists a unique solution given by

$$T = \sum_{i < 0} e^{\lambda_i z} (B_- \phi_{\mathcal{D}} | \phi_i)_{\mathcal{H}} U_i.$$

In this case the temperature at infinity is 0.

b) If the constants are not controlled and $\int_{\Omega} h \neq 0$, then there exists a unique solution given by

$$T = \sum_{i < 0} e^{\lambda_i z} (B_- (\phi_{\mathcal{D}} - T_{\infty} \Phi) | \phi_i)_{\mathcal{H}} U_i + T_{\infty},$$

where T_{∞} is an arbitrary constant in the case $\int_{\Omega} h < 0$ and $T_{\infty} = (\Phi - PB_- \Phi | \phi_{\mathcal{D}})_{\mathcal{H}} (\Phi - PB_- \Phi | \Phi)_{\mathcal{H}}^{-1}$ in the case $\int_{\Omega} h > 0$. In this case the temperature at infinity is the constant T_{∞} . Note that if $\int_{\Omega} h > 0$, the temperature at infinity is determined by $\phi_{\mathcal{D}}$ whereas in the case $\int_{\Omega} h < 0$, it is a free parameter of the problem.

c) In the balanced case, the set of solutions is given by:

$$T(z) = \sum_{i < 0} (B_- (\phi_{\mathcal{D}} - c_1 \Phi - c_2 (\Delta_{\sigma}^{-1} h, 0)) | \phi_i)_{\mathcal{H}} U_i e^{\lambda_i z} + c_1 + c_2 (z + \Delta_{\sigma}^{-1} h),$$

where c_2 is an arbitrary constant and

$$c_1 = (\Phi - PB_- \Phi | \phi_{\mathcal{D}} - c_2 (\Delta_{\sigma}^{-1} h, 0))_{\mathcal{H}} (\Phi - PB_- \Phi | \Phi)_{\mathcal{H}}^{-1}.$$

In this case the temperature at infinity has the linear growth rate $T(z) \simeq c_2 z + (c_2 \Delta_{\sigma}^{-1} h + c_1) + o(z)$. If the temperature is not allowed to have a linear growth rate, then $c_2 = 0$ and the temperature at infinity is c_1 which is determined by the initial conditions.

Proof. We use the result of Section 2.2 on the solution of the evolution equation, that is T solves (E) if and only if there exists $\psi \in \mathcal{R}(\mathcal{A})$ and constants c_1 and c_2 such that:

$$T(z) = \sum_{i \in \mathbb{Z}^*} (\psi | \phi_i)_{\mathcal{H}} U_i e^{\lambda_i z} + c_1 + c_2 (z + \Delta_{\sigma}^{-1} h),$$

where $c_1 = c_2 = 0$ if the constants are controlled and $c_2 = 0$ in the non-balanced case. The subexponential growth condition ensures that $\pi_+ \psi = 0$. The condition $(T(z=0), 0) = \phi_{\mathcal{D}}$ yields

$$P\phi_{\mathcal{D}} = P\psi + c_1 \Phi + c_2 (\Delta_{\sigma}^{-1} h, 0). \quad (4.2)$$

Using Theorem 5 leads to distinguishing the following cases:

- a) If the constants are controlled then $c_1 = c_2 = 0$ and the equation $P\phi_{\mathcal{D}} = P\psi$ with $\pi_+ \psi = 0$ has the unique solution $\psi = B_- \phi_{\mathcal{D}}$.
- b) In the non-balanced case, $c_2 = 0$ and $P(\phi_{\mathcal{D}} - c_1 \Phi) = P\psi$ together with $\pi_+ \psi = 0$ implies $\psi = B_- (\phi_{\mathcal{D}} - c_1 \Phi)$ without any additional assumption in the case $\int_{\Omega} h < 0$. In the case where $\int_{\Omega} h > 0$ the compatibility condition is

$$(\Phi - PB_- \Phi | \phi_{\mathcal{D}} - c_1 \Phi) = 0, \text{ which gives } c_1 = (\Phi - PB_- \Phi | \phi_{\mathcal{D}}) (\Phi - PB_- \Phi | \Phi)^{-1}.$$

c) Finally, in the balanced case, let us fix an arbitrary value c_2 . The conditions $\psi \in \mathcal{R}(\mathcal{A})$ and $\pi_+\psi = 0$ are equivalent to $\psi \in \mathcal{R}(\pi_-)$. Denoting $\tilde{\Phi} = (\Delta_\sigma^{-1}h, 0)$ we re-cast (4.2) into

$$P\phi_{\mathcal{D}} = P\psi + c_1\Phi + c_2\tilde{\Phi}, \quad \psi \in \mathcal{R}(\pi_-)$$

or equivalently

$$P(\phi_{\mathcal{D}} - c_2\tilde{\Phi}) = P(\psi + c_1\phi_0), \quad \psi \in \mathcal{R}(\pi_-).$$

In view of Theorem 5 with $\phi = \phi_{\mathcal{D}} - c_2\tilde{\Phi}$, there is a unique solution to the above equation given by

$$c_1 = (\Phi - PB_- \Phi | \phi_{\mathcal{D}} - c_2\tilde{\Phi})_{\mathcal{H}} (\Phi - PB_- \Phi | \Phi)_{\mathcal{H}}^{-1}, \quad \psi = B_-(\phi_{\mathcal{D}} - c_1\Phi - c_2\tilde{\Phi}).$$

177 \square

4.2. Semi-infinite problem, Neumann Inlet/Outlet Condition. We consider the Neumann (I/OBC) condition:

$$\partial_z T|_{z=0} = S_0. \quad (4.3)$$

178 Denote $\phi_{\mathcal{N}} = (S_0, 0) \in \mathcal{H}$.

179 **PROPOSITION 7.** *Consider the Graetz problem (E) on the semi-infinite cylinder*
 180 $\Omega \times [0, +\infty)$, *with subexponential growth together with Neumann Inlet/Outlet condition*
 181 *(4.3).*

a) *If the constants are controlled, then there exists a unique solution given by*

$$T = \sum_{i < 0} e^{\lambda_i z} (\mathcal{A}^{-1} B_- \phi_{\mathcal{N}} | \phi_i)_{\mathcal{H}} U_i.$$

182 *In this case the temperature at infinity is 0.*

b) *If the constants are not controlled and $\int_{\Omega} h \neq 0$. If $\int_{\Omega} h > 0$ there always exists a solution, if $\int_{\Omega} h < 0$, there exists a solution if and only if*

$$(\Phi - PB_- \Phi | \phi_{\mathcal{N}}) = 0.$$

When the solution exists, it is of the form

$$T = \sum_{i < 0} e^{\lambda_i z} (\mathcal{A}^{-1} B_- (\phi_{\mathcal{N}} - T_{\infty} \Phi) | \phi_i)_{\mathcal{H}} U_i + T_{\infty},$$

183 *where the temperature at infinity T_{∞} is a free parameter of the problem.*

c) *In the balanced case, the set of solutions is given by*

$$T(z) = \sum_{i < 0} (\mathcal{A}^{-1} B_- (\phi_{\mathcal{N}} + c_2 \Phi) | \phi_i)_{\mathcal{H}} U_i e^{\lambda_i z} + c_1 + c_2(z + \Delta_\sigma^{-1}h),$$

where c_1 is an arbitrary constant and c_2 is given by:

$$c_2 = -(\Phi - PB_- \Phi | \phi_{\mathcal{N}}) (\Phi - PB_- \Phi | \Phi)^{-1}.$$

184 *In this case the temperature at infinity has the linear growth rate $T(z) \simeq$*
 185 $c_2 z + (c_2 \Delta_\sigma^{-1}h + c_1) + o(z)$.

Proof. We proceed as in the previous section. It follows from the result of Section 2.2 on the solution of the evolution equation that T solves (E) if and only if there exists $\psi \in \mathcal{R}(\mathcal{A})$ and constants c_1 and c_2 such that:

$$T(z) = \sum_{i \in \mathbb{Z}^*} (\psi | \phi_i)_{\mathcal{H}} U_i e^{\lambda_i z} + c_1 + c_2(z + \Delta_{\sigma}^{-1} h),$$

where $c_1 = c_2 = 0$ if the constants are controlled and $c_2 = 0$ in the non-balanced case. Derivating w.r.t. z one finds

$$P\phi_{\mathcal{N}} = P\mathcal{A}\psi + c_2\Phi, \quad \pi_+\psi = 0. \quad (4.4)$$

Using Theorem 5 leads to distinguishing the following cases:

- a) If the constants are controlled then $c_1 = c_2 = 0$ and equation (4.4) admits a unique solution $\mathcal{A}\psi = B_-\phi_{\mathcal{N}}$. The invertibility of \mathcal{A} gives the result.
- b) If the constants are not controlled then $c_2 = 0$. If $\int_{\Omega} h > 0$ there is always a solution $\mathcal{A}\psi$ to equation (4.4) and the operator \mathcal{A} is invertible, hence there exists a unique solution ψ to (4.4) given by $\psi = \mathcal{A}^{-1}B_-\phi_{\mathcal{N}}$. The constant c_1 is then a free parameter of the problem. If $\int_{\Omega} h < 0$ then the condition for equation (4.4) to admit a solution is:

$$(\Phi - PB_-\Phi | \phi_{\mathcal{N}}) = 0.$$

If this condition is met, by the invertibility of \mathcal{A} , $\psi = \mathcal{A}^{-1}B_-\phi_{\mathcal{N}}$ is the unique solution to (4.4) and c_1 is a free parameter of the problem.

- c) In the balanced case let c_2 be an arbitrary constant. It follows from Theorem 5 that $\mathcal{A}\psi$ satisfies (4.4) if and only if

$$\mathcal{A}\psi = B_-P\phi + \frac{(\Phi - PB_-\Phi | \phi)_{\mathcal{H}}}{(\Phi - PB_-\Phi | \Phi)_{\mathcal{H}}} (B_-\Phi - \phi_0)_{\mathcal{H}} \quad \text{where } \phi = \phi_{\mathcal{N}} + c_2\Phi. \quad (4.5)$$

For ψ to exist, the right hand side must belong to the range of \mathcal{A} , i.e. be orthogonal to ϕ_0 . Performing the scalar product of the left hand side of (4.5) with ϕ_0 and recalling that the range of B_- is orthogonal to $\mathcal{K}(\mathcal{A})$ we obtain the following necessary condition:

$$(\Phi - PB_-\Phi | \phi_{\mathcal{N}} + c_2\Phi)_{\mathcal{H}} = 0,$$

which is equivalent to:

$$c_2 = -(\Phi - PB_-\Phi | \phi_{\mathcal{N}})(\Phi - PB_-\Phi | \Phi)^{-1}.$$

186 Conversely, if the above condition is met, then equation (4.5) admits a unique inverse
187 in $\mathcal{R}(\mathcal{A})$ and c_1 is a free parameter of the problem. \square

5. Resolution of the problem in a finite domain. We aim to solve Graetz equation in a domain of finite length $\Omega \times [-L, L]$:

$$\begin{cases} c\partial_{zz}T + \operatorname{div} \sigma \nabla T - h\partial_z T = 0 & \Omega \times [-L, L], \\ \text{(LBC)} & (\partial\Omega) \times [-L, L], \\ \text{Inlet/Outlet condition} & \Omega \times \{-L, L\}, \end{cases} \quad (5.1)$$

where the I/O conditions can be of Neumann or Dirichlet type. According to Section 2.2, the solutions may be sought in the form

$$T(z) = \sum_{i < 0} (\psi | \phi_i) e^{\lambda_i(z+L)} U_i + \sum_{i > 0} (\psi | \phi_i) e^{\lambda_i(z-L)} U_i + c_1 + c_2(z + \Delta_{\sigma}^{-1} h),$$

with $c_1 = c_2 = 0$ if the constants are controlled and $c_2 = 0$ in the non-balanced case.

The unknowns in this equation are $(\psi|\phi_i)$ for $i < 0$ and $i > 0$, plus possibly (depending on the case) c_1 and c_2 . Note that $\sum_{i<0}(\psi|\phi_i)U_i = P\pi_-\psi$, and therefore if $P\pi_-\psi$ is known it suffices to decompose this vector on the basis of $L^2(\Omega)$ given by $(U_i)_{i<0}$ to obtain the desired coefficients for $i < 0$. Similarly the coefficients $(\psi|\phi_i)$ for $i > 0$ are obtained by considering the coefficients of $P\pi_+\psi$ on the basis composed of the $(U_i)_{i>0}$. Therefore the unknowns to be determined are $P\pi_-\psi, P\pi_+\psi$ plus possibly c_1 and c_2 .

Let X be the vector composed of all the unknowns. Then satisfying the I/O conditions amounts to solving a linear system for X . In the rest of this section we detail the linear system in each case, but beforehand we focus on a linear operator involved in the system.

5.1. Study of the linear operator M . We define and study a linear operator that will be involved in the solution of the problem in a cylinder of length $2L$.

PROPOSITION 8. *Let M_\pm be the operators from $\mathcal{R}(P)$ to $\mathcal{R}(P)$ and M be given by:*

$$M_\pm = Pe^{\mp 2LA}B_\pm \quad \text{and} \quad M = \begin{pmatrix} 0 & M_+ \\ M_- & 0 \end{pmatrix}.$$

Then

a) *there exists a constant C such that*

$$\|M\| \leq Ce^{-2\lambda L}, \quad \text{where } \lambda = \min(\lambda_1, -\lambda_{-1}).$$

As a consequence $\|M\| < 1$ for sufficiently large L .

b) *If the constants are controlled, then for L positive sufficiently small, $\|M^2\| < 1$.*

c) *It follows that $Id + M$ is invertible on $\mathcal{R}(P) \times \mathcal{R}(P)$ for large L and for small positive L .*

Proof.

a) Since $M_+ = Pe^{-2LA}B_+$ we have

$$\|M_+\| \leq \|B_+\|e^{-2L\lambda_1}.$$

A similar upper bound for M_- gives the result.

b) Define

$$J(L) = \sup_{\|(\phi_1, \phi_2)\|=1} \|M^2(\phi_1, \phi_2)\| < 1$$

Since $J(0) = 1$, it is sufficient to prove that $J'(0) < 0$. Note that

$$M^2 = \begin{pmatrix} M_+M_- & 0 \\ 0 & M_-M_+ \end{pmatrix}.$$

Let us fix $\phi \in \mathcal{R}(P)$, and define $j(L) = \|M_+(L)M_-(L)\phi\|^2$. Then $j(0) = \|\phi\|^2$ and it remains to prove that $j'(0) \leq -C\|\phi\|^2$ with a positive constant C independent of ϕ .

The derivative of j is:

$$\begin{aligned} j'(L) &= (M'_+(L)M_-(L)\phi + M_+(L)M'_-(L)\phi|M_+(L)M_-(L)\phi) \\ &= (-2PAe^{-2LA}B_+Pe^{2LA}B_-\phi + 2Pe^{-2LA}B_+PAe^{2LA}B_-\phi|Pe^{-2LA}B_+Pe^{2LA}B_-\phi) \end{aligned}$$

213 hence

$$j'(0) = -2(PAB_+\phi|\phi) + 2(PAB_-\phi|\phi).$$

214 But since $P\phi = \phi$, $PB_+P = P$ and $P\mathcal{A} = \mathcal{A} + PAP - \mathcal{A}P$ we have

$$\begin{aligned} (PAB_+\phi|\phi) &= (PAB_+\phi|B_+\phi) \\ &= ((\mathcal{A} + PAP - \mathcal{A}P)B_+\phi|B_+\phi) \\ &= (\mathcal{A}B_+\phi|B_+\phi) + (PAPB_+\phi|B_+\phi) - (\mathcal{A}PB_+\phi|B_+\phi) \\ &= (\mathcal{A}B_+\phi|B_+\phi) + (\mathcal{A}\phi|\phi) - (\mathcal{A}\phi|B_+\phi). \end{aligned}$$

This proves that

$$(PAB_+\phi|\phi) = (\mathcal{A}\phi|B_+\phi) = \frac{1}{2}((\mathcal{A}B_+\phi|B_+\phi) + (\mathcal{A}\phi|\phi)).$$

Similarly we obtain that

$$(PAB_-\phi|\phi) = \frac{1}{2}((\mathcal{A}B_-\phi|B_-\phi) + (\mathcal{A}\phi|\phi)).$$

As a summary we find that

$$j'(0) = -(\mathcal{A}B_+\phi|B_+\phi) + (\mathcal{A}B_-\phi|B_-\phi) < \lambda_{-1}\|B_-\phi\|^2 - \lambda_1\|B_+\phi\|^2 < (\lambda_{-1} - \lambda_1)\|\phi\|^2.$$

215 c) The operator $Id + M$ is invertible for large L by a). Note that M_\pm as endomorphism
 216 of $\mathcal{R}(P)$ are compact for $L > 0$ and equal to identity for $L = 0$. As a result $Id + M$
 217 is invertible for small $L > 0$ if and only if there is no eigenvector associated to the
 218 value -1 . A sufficient condition for invertibility is then that M^2 does not admit 1 as
 219 eigenvalue, which is proved in b) for L sufficiently small. \square

220 **5.2. The Dirichlet case.** The different cases for Dirichlet I/O condition are
 221 summarized in the following

PROPOSITION 9. *The Dirichlet I/O condition $T|_{z=-L} = T_{-L}$ and $T|_{z=L} = T_{+L}$ are equivalent to the following linear system*

$$ZX = b,$$

222 where Z , X and b are defined depending of the (LBC) and given in the table below.

	constants controlled	constants not controlled	
		unbalanced ($\int_\Omega h > 0$)	balanced
Z	$\begin{pmatrix} Id & M_+ \\ M_- & Id \end{pmatrix}$	$\begin{pmatrix} Id & M_+ & \Phi \\ M_- & Id & \Phi \\ 0 & M_+^* u_-^T & (\Phi u_-) \end{pmatrix}$	$\begin{pmatrix} Id & M_+ & \Phi & -L\Phi + \tilde{\Phi} \\ M_- & Id & \Phi & L\Phi + \tilde{\Phi} \\ 0 & M_+^* u_-^T & (\Phi u_-) & (-L\Phi + \tilde{\Phi} u_-) \\ M_-^* u_+^T & 0 & (\Phi u_+) & (L\Phi + \tilde{\Phi} u_+) \end{pmatrix}$
X	$\begin{pmatrix} P\pi_- \psi \\ P\pi_+ \psi \end{pmatrix}$	$\begin{pmatrix} P\pi_- \psi \\ P\pi_+ \psi \\ c_1 \end{pmatrix}$	$\begin{pmatrix} P\pi_- \psi \\ P\pi_+ \psi \\ c_1 \\ c_2 \end{pmatrix}$
b	$\begin{pmatrix} \phi_{-L} \\ \phi_{+L} \end{pmatrix}$	$\begin{pmatrix} \phi_{-L} \\ \phi_{+L} \\ (\phi_{-L} u_-) \end{pmatrix}$	$\begin{pmatrix} \phi_{-L} \\ \phi_{+L} \\ (\phi_{-L} u_-) \\ (\phi_L u_+) \end{pmatrix}$

where we recall that $\tilde{\Phi} = (\Delta_\sigma^{-1}h, 0)$ and we define $\phi_{\pm L} = (T_\pm, 0)$, $u_- = \Phi - PB_- \Phi$ and $u_+ = \Phi - PB_+ \Phi$.

Moreover, for sufficiently large L this system is invertible.

Note 1: Thanks to the Lax-Milgram theorem in 3D (see Section 1.3), we know before hand that there exists a unique solution to the system $ZX = b$.

Note 2: In the case when the constants are not controlled and $\int_\Omega h < 0$ it suffices to change the sign of z , or equivalently to remplace the $-$ by $+$.

Proof. The (I/OBC) are equivalent to the following

$$\begin{cases} P\pi_- \psi = P\theta_- & \text{with } \theta_- = \phi_{-L} - e^{-2LA}\pi_+ \psi - c_1 \Phi - c_2(-L\Phi + \tilde{\Phi}), \\ P\pi_+ \psi = P\theta_+ & \text{with } \theta_+ = \phi_{+L} - e^{2LA}\pi_- \psi - c_1 \Phi - c_2(L\Phi + \tilde{\Phi}). \end{cases} \quad (5.2)$$

Combining

$$Pe^{-2LA}\pi_+ \psi = Pe^{-2LA}B_+ P\pi_+ \psi = M_+ P\pi_+ \psi,$$

and the similar version when the roles of $+$ and $-$ are interchanged with equations (5.2) we obtain the first two rows of the matrix Z .

a) when the constants are controlled, $c_1 = c_2 = 0$ and (5.2) reads $ZX = b$.

b) when the constants are not controlled and $\int_\Omega h > 0$ then $c_2 = 0$. Theorem 5 requires an additional compatibility condition to solve the first equation. This condition reads

$$(\Phi - PB_- \Phi | \theta_-) = 0,$$

which is the additional equation in the system $ZX = b$.

c) In the balanced case, after the change of variable $\tilde{\psi} = \pi_- \psi$ the first equation in (5.2) $P\pi_- \psi = P\theta_-$ is equivalent to

$$\begin{cases} P\tilde{\psi} = P\theta_- \\ \pi_+ \tilde{\psi} = 0 \end{cases} \quad \text{and} \quad (\tilde{\psi} | \phi_0) = 0.$$

Theorem 5 gives an explicit expression for the solution of the system on the left, and the condition on the right becomes $(\theta_- | \Phi - PB_- \Phi) = 0$, which is the third row of the system $ZX = b$. The last row is obtained using the second equation in (5.2).

When L becomes large, M_+ and M_- are exponentially small and in each case the matrix Z is asymptotic to an invertible matrix. The sole non-obvious case is the balanced case, where one can observe that the 2×2 lower right block is asymptotically

equivalent to $\begin{pmatrix} (\Phi | u_-) & (-L\Phi | u_-) \\ (\Phi | u_+) & (L\Phi | u_+) \end{pmatrix}$ which has a determinant $2L(\Phi | u_-)(\Phi | u_+) \neq 0$.

When L is large, Z can be rewritten as $Z = A + B$ with B small and A easily inverted. One can use a Neumann series strategy to solve $Zx = b$. \square

5.3. The Neumann Inlet/Outlet case. The different cases for Neumann I/O condition are summarized in the following

PROPOSITION 10. *The Neumann I/O condition $\partial_z T|_{z=-L} = S_{-L}$ and $\partial_z T|_{z=L} = S_{+L}$ are equivalent to the following linear system*

$$ZX = b,$$

where Z , X and b are defined depending of the (LBC) and given in the table below.

	constants controlled	constants not controlled	
		unbalanced ($\int_{\Omega} h > 0$)	balanced
Z	$\begin{pmatrix} Id & M_+ \\ M_- & Id \end{pmatrix}$	$\begin{pmatrix} Id & M_+ \\ M_- & Id \\ 0 & M_+^* u_-^T \end{pmatrix}$	$\begin{pmatrix} Id & M_+ & \Phi \\ M_- & Id & \Phi \\ 0 & M_+^* u_-^T & (\Phi u_-) \\ M_-^* u_+^T & 0 & (\Phi u_+) \end{pmatrix}$
X	$\begin{pmatrix} P\mathcal{A}\pi_- \psi \\ P\mathcal{A}\pi_+ \psi \end{pmatrix}$	$\begin{pmatrix} P\mathcal{A}\pi_- \psi \\ P\mathcal{A}\pi_+ \psi \end{pmatrix}$	$\begin{pmatrix} P\mathcal{A}\pi_- \psi \\ P\mathcal{A}\pi_+ \psi \\ c_2 \end{pmatrix}$
b	$\begin{pmatrix} \phi_{-L} \\ \phi_{+L} \end{pmatrix}$	$\begin{pmatrix} \phi_{-L} \\ \phi_{+L} \\ (\phi_{-L} u_-) \end{pmatrix}$	$\begin{pmatrix} \phi_{-L} \\ \phi_{+L} \\ (\phi_{-L} u_-) \\ (\phi_L u_+) \end{pmatrix}$

where we define $\phi_{\pm L} = (S_{\pm}, 0)$, $u_- = \Phi - PB_- \Phi$ and $u_+ = \Phi - PB_+ \Phi$.

Note 1: when the constants are not controlled the value of c_1 is arbitrary. In these cases the linear systems are rectangular and the existence of the solution depends on a compatibility conditions that expresses that b is in the range of Z .

Note 2: once the quantities $P\mathcal{A}\pi_{\pm}\psi$ are known, then the $(\mathcal{A}\psi|\phi_i)$ for $i > 0$ and $i < 0$ can be computed as explained above, and $(\psi|\phi_i)$ is obtained by dividing by λ_i .

Proof. The (I/OBC) are equivalent to the following

$$\begin{cases} P\mathcal{A}\pi_- \psi = P\theta_- & \text{with } \theta_- = \phi_{-L} - e^{-2L\mathcal{A}}\mathcal{A}\pi_+ \psi - c_1\Phi - c_2(-L\Phi + \tilde{\Phi}), \\ P\mathcal{A}\pi_+ \psi = P\theta_+ & \text{with } \theta_+ = \phi_{+L} - e^{2L\mathcal{A}}\mathcal{A}\pi_- \psi - c_1\Phi - c_2(L\Phi + \tilde{\Phi}). \end{cases} \quad (5.3)$$

A discussion similar to the Dirichlet case leads to the result. \square

6. Numerical tests.

6.1. First test case: a domain of finite length. The section of the domain of the first test-case is the square $\Omega = [-5, 5]^2$ with a circular fluid subdomain of radius 2 centered at the origin. The velocity and eigenvalues of the operator \mathcal{A} are computed with $P1$ finite element methods on the mesh of Figure 6.1 . The velocity has

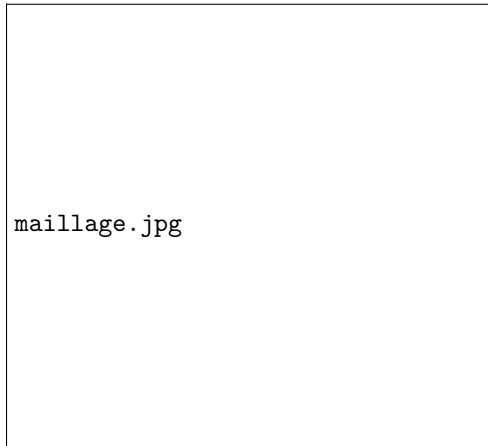


FIGURE 6.1. The mesh for the first test case is composed of 13589 vertices and 26776 triangles. The solid domain is in white and the fluid domain in grey.

a parabolic profile (Poiseuille flow) with prescribed total flow $Q \in \{1, 10, 100, 1000\}$. The lateral boundary conditions are of Robin type with parameter a . The thermal conductivities are equal to $c = \sigma = 1$. In total 100 eigenvalues / eigenvectors of \mathcal{A} are computed.

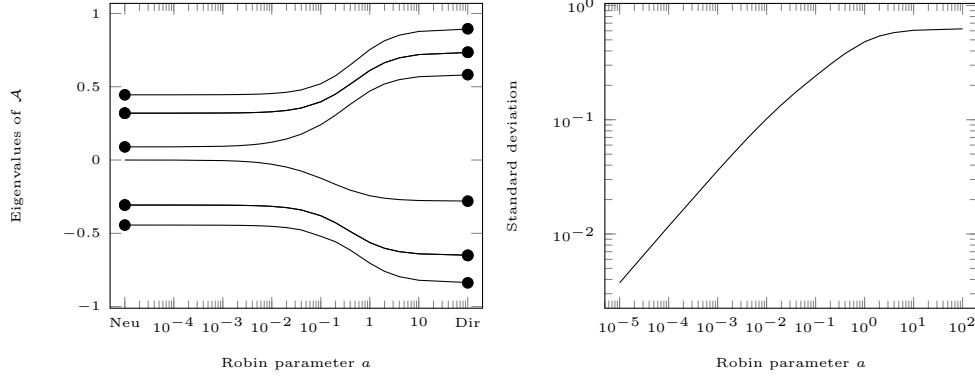


FIGURE 6.2. On the left: Evolution of the eigenvalues of \mathcal{A} of smallest magnitude for varying parameter of the Robin lateral boundary condition. Eigenvalues for the Neumann (resp. Dirichlet) boundary conditions are shown as bullet on the left (resp. right) of the curves. On the right: relative L^2 difference between the eigenvector with largest negative eigenvalue and its mean as the Robin parameter varies.

We first set $Q = 10$ and vary the Robin parameter a . When $a = 0$, one retrieves the Neumann case and when $a = +\infty$, one retrieves the Dirichlet case. In order to emphasize this fact we plot in Figure 6.2 (left) the eigenvalues of smallest magnitude for different values of a . We also plot with dots the eigenvalues associated to the Neumann problem (on the left of the curves) and the one associated to the Dirichlet case (on the right of the curves). The smooth transition from Neumann to Dirichlet as the Robin parameter varies is striking except from the fact that there exists an eigenvalue that goes to zero as a goes to zero even if the Neumann problem does not have zero as eigenvalue. We claim that this behavior is consistent with theory. First 0 is not an eigenvalue of the Neumann case since the total flow is non-zero (hence we are not in a balanced case even if the constants are not controlled). Second, we remark that the zero eigenvalue is the limit of a negative eigenvalue. Remember from Proposition 6 that it is always possible to decompose a temperature field on the set of negative eigenvectors in the Robin case (part a), but for the Neumann case it is necessary to add a constant (part b). In other words, in the Neumann case the constant must be added to the negative eigenvectors to obtain a Hilbert basis of H , while the set of positive eigenvectors form a Hilbert basis on their own. This explains why the constant emerges as the limit of a negative eigenvector, see Figure 6.2 (right) where the convergence of the eigenvector to the constant is numerically demonstrated.

In a second parametric study, we fix $a = 1$ and we let both Q and L vary. First we plot the spectral radius of the matrix M defined in Proposition 8 versus the exchanger length L for the different values of the total flow Q in Figure 6.3 (left). Figure 6.3 (center) shows the evolution of the 5 smallest positive and 5 largest negative eigenvalues of M for a fixed total flow $Q = 20$. This test-case shows that, apart from the case $L = 0$, the spectral radius of the matrix M is always smaller than

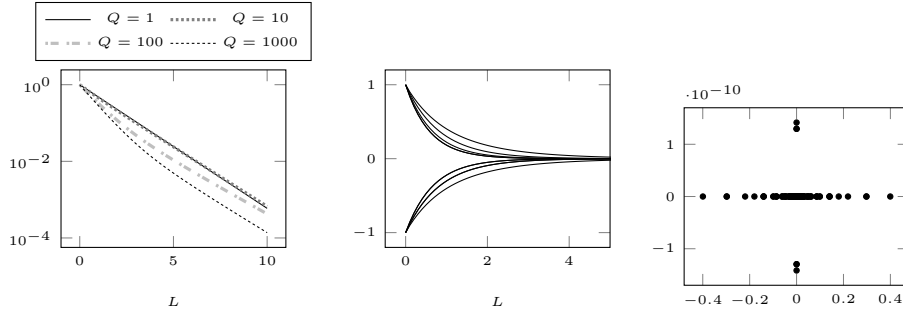


FIGURE 6.3. Eigenvalues of M for a Robin test case with $a = 1$. Left: evolution of the spectral radius for different values of the total flow Q . Center: the evolution of the five largest positive and first smallest negative eigenvalue for a total flow of $Q = 20$. Right: eigenvalues in the complex plane for $L = 1$ and $Q = 20$.

one, so that the matrix $Id + M$ is indeed always invertible. The exponential decrease for large L and the decrease at the origin follows from Proposition 8. Moreover, since the spectral radius of M is strictly smaller than one, a Neumann series strategy to solve

$$(Id + M)^{-1}b = \sum_k (-M)^k b$$

is legitimate. In Figure 6.3 (right), the whole spectrum of M is shown in the complex plane. Although the spectrum seems real, we do not have mathematical proof of this fact.

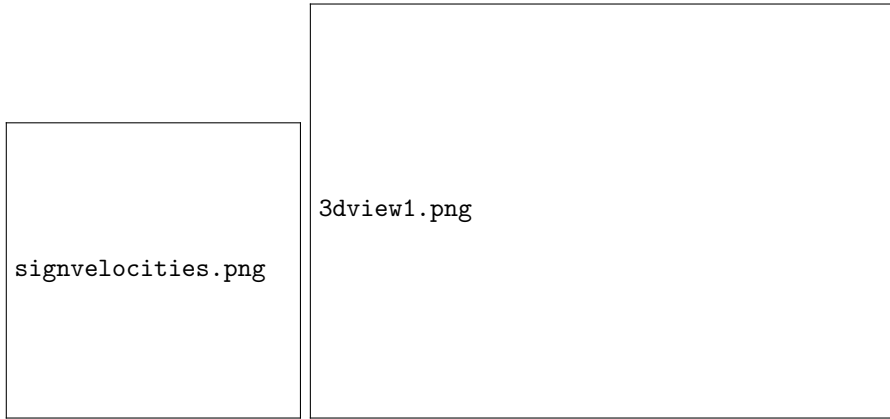


FIGURE 6.4. Geometry of the periodic exchanger. On the left, a cut inside the exchanger with the sign of the fluid velocities. On the right, a 3D representation of the exchanger. The tubes where the temperature is set at ∞ are colored accordingly to their temperature.

6.2. Second test-case: a periodic exchanger. The second test case consists of a heat exchanger with periodic boundary conditions. The whole device consists of one solid exchanger through which pass four tubes containing fluids. A cut along the middle of the exchanger is shown in Figure 6.4 (left) where the sign of the velocity of the fluid in the inner tubes is displayed. The fluids are assumed to obey a Poiseuille



FIGURE 6.5. *Four different solutions of the periodic exchanger, with different length and total flow. The length of the exchanger is set to $L = 10$ on top and $L = 20$ on bottom. The total flow is set to $Q = 10$ on left and $Q = 30$ on right.*

flow, the velocities are then quadratic in the radial coordinates of their corresponding tubes. The length of the exchanger is denoted L , the section of the exchanger is the square $[-4, 4]^2$, the radii of the inner tubes are fixed to 1 and the distance of the center of the inner tubes to the center of the exchanger is $\sqrt{2^2 + 2^2}$. The conductivity in both the fluid and solid part is set to 1. The temperature is fixed for the four tubes with incoming flow (two at each side) on the exchanger, the warm temperature is set to $+1$ and the cold temperature to -1 , see Figure 6.4 (left). In what follows, Q denotes the total flow of fluid in one tube. We glue together the different Graetz problems using the methodology developped in [4].

In Figure 6.5, four solutions are shown for different values of the length L and the flow Q .

Figure 6.6 displays the efficiency and the total exchange for different values of Q and L . For a tube containing fluid whose velocity is positive (resp. negative), the temperature at $-\infty$ (resp. $+\infty$) is set to 1 (resp -1), the efficiency of the exchanger is then defined by $-T_{+\infty}/T_{-\infty}$ (resp $-T_{-\infty}/T_{+\infty}$), where $T_{\pm\infty}$ is the temperature at infinity. This efficiency is between -1 and 1. The exchange is simply the total amount of heat exchanged and is equal to Q times the efficiency. The aim of this test-case is

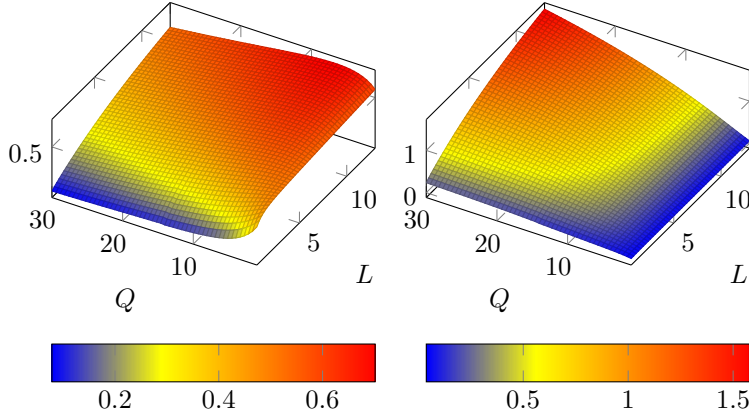


FIGURE 6.6. Values of the efficiency (left) and the exchange (right) for different values of the length L and the total flow Q of the periodic exchanger. L ranges from 0.5 to 13 and Q ranges from 1 to 30. Each direction has been sampled 50 times for a total of 2500 exchanger computations.

to document the fact that our method is able to deal with any boundary conditions and type of exchanger. It is well suited for parametric studies.

7. Conclusion. In the present work we have proposed a general framework dedicated to the resolution of the generalized Graetz problem in arbitrary geometry, involving any type of boundary conditions. The main novelty is the introduction of insulating boundary condition (Neumann or periodic) that allow to model realistic heat exchangers. Our study highlighted a special case that we call the balanced case, when $\int_{\Omega} h = 0$ (together with Neumann or periodic boundary condition) where the solution is different than in the general case. We have also proposed a number of numerical illustration in various test cases.

Appendix A. Technical lemmas. We prove here results that will be used in the sequel.

LEMMA 11. For each $\phi = (u, s) \in \mathcal{D}(\mathcal{A})$, $\tilde{\phi} = (\tilde{u}, \tilde{s}) \in \mathcal{H}$, we have

$$(Id - P)\mathcal{A}(Id - P)\phi = 0. \quad (\text{A.1})$$

$$(P\mathcal{A}P\phi|\tilde{\phi})_{\mathcal{H}} = \int_{\Omega} hu\tilde{u} \leq \|h\|_{L^{\infty}(\Omega)} \|P\phi\|_{\mathcal{H}} \|P\tilde{\phi}\|_{\mathcal{H}}. \quad (\text{A.2})$$

Proof. This results from elementary calculations using the definition of \mathcal{A} (2.1), and the definition of P (3.1). \square

LEMMA 12. Let $\Phi = (1, 0) \in \mathcal{H}$. Let $\phi \in \mathcal{D}(\mathcal{A}^{-1/2})$ such that $P\phi = \phi$, then

$$(\mathcal{A}^{-1}\phi|\phi)_{\mathcal{H}} = \left(\int_{\Omega} h \right)^{-1} (\phi|\Phi)_{\mathcal{H}}^2 \text{ in the non-balanced case}$$

$$(\mathcal{A}^{-1}\phi|\phi)_{\mathcal{H}} = 0 \text{ in the "balanced" or "constant controlled" case}$$

Proof. The expression of \mathcal{A}^{-1} is given in section 2.1 for the various cases. Let $\phi \in \mathcal{D}(\mathcal{A}^{-1/2})$ such that $P\phi = \phi$, hence there exists $u \in L^2(\Omega)$ such that $\phi = (u, 0)$. If the constants are controlled then $\mathcal{A}^{-1}\phi = (0, \Delta_{\sigma}^{-1}(-cu))$ and $(\mathcal{A}^{-1}\phi|\phi)_{\mathcal{H}} = 0$. If

the constants are not controlled and $\int_{\Omega} h \neq 0$, then

$$\mathcal{A}^{-1}\phi = (k, \Delta_{\sigma}^{-1}(-cu + hk)) \text{ with } k \int_{\Omega} h = \int_{\Omega} cu,$$

hence

$$(\mathcal{A}^{-1}\phi|\phi)_{\mathcal{H}} = \int_{\Omega} cku = (\int_{\Omega} h)^{-1}(\int_{\Omega} cu)^2.$$

326 Finally, in the balanced case, since $\phi \in \mathcal{D}(\mathcal{A}^{-1/2})$, then $(\phi|\phi_0)_{\mathcal{H}} = 0$ and $\mathcal{A}^{-1}\phi =$
 327 $(0, \Delta_{\sigma}^{-1}(-cu)) + k\phi_0$ exists and $(\mathcal{A}^{-1}\phi|\phi)_{\mathcal{H}} = 0$. \square

LEMMA 13. *In the balanced case,*

$$(\Phi - PB_{-}\Phi|\Phi)_{\mathcal{H}} \neq 0.$$

328

329 *Proof.* Suppose the contrary and set $\theta = \Phi - PB_{-}\Phi$, we have $P\theta = \theta$ and by
 330 definition of B_{-} , we have $\pi_{-}\theta = 0$. Moreover, we have

$$\begin{aligned} (\theta|\phi_0)_{\mathcal{H}} &= (\Phi - PB_{-}\Phi|\phi_0)_{\mathcal{H}} = (P\Phi - PB_{-}\Phi|\phi_0)_{\mathcal{H}} = (\Phi - B_{-}\Phi|P\phi_0)_{\mathcal{H}} \\ &= (\Phi - B_{-}\Phi|P\Phi)_{\mathcal{H}} = (\Phi - PB_{-}\Phi|\Phi)_{\mathcal{H}} = 0 \end{aligned}$$

Lemma 12 ensures that $\mathcal{A}^{-1}\theta$ exists and that

$$(\mathcal{A}^{-1}\theta|\theta)_{\mathcal{H}} = 0$$

Since θ belongs to $\mathcal{R}(\pi_{+})$ and all the eigenvalues of \mathcal{A} are positive on this space, this implies that $\theta = 0$ and then $\Phi = PB_{-}\Phi$. Hence there exists s such that $B_{-}\Phi = (1, s)$ and

$$(\mathcal{A}B_{-}\Phi|B_{-}\Phi)_{\mathcal{H}} = ((h - \Delta_{\sigma}s, 0) | (1, s))_{\mathcal{H}} = 0.$$

331 But $B_{-}\Phi$ belongs to $\mathcal{R}(\pi_{-})$ and since all the eigenvalues of \mathcal{A} are negative on $\mathcal{R}(\pi_{-})$,
 332 $(\mathcal{A}B_{-}\Phi|B_{-}\Phi)_{\mathcal{H}} = 0$ implies that $B_{-}\Phi = 0$ which is in violation of $\Phi = PB_{-}\Phi$. Hence
 333 $(\Phi - B_{-}\Phi|\Phi)_{\mathcal{H}} \neq 0$. \square

Appendix B. Proof of Theorem 4. Let $M \in \mathbb{N}^*$ and denote for short $\pi = \pi_{[-M, -1]}$. The operator $\pi P \pi$ is a symmetric operator on a finite-dimensional space, hence it is diagonalisable in an orthonormal basis. The first step is to prove that this operator is definite positive with a lower bound on its eigenvalues that is independent of M . Let ρ be an eigenvalue of $\pi P \pi$ and \mathbf{v} an associated normalized eigenvector: $\pi P \pi \mathbf{v} = \rho \mathbf{v}$, $(\mathbf{v}|\mathbf{v})_{\mathcal{H}} = 1$ and $\pi \mathbf{v} = \mathbf{v}$. Since

$$\rho = (\pi P \pi \mathbf{v}|\mathbf{v})_{\mathcal{H}} = (P \pi \mathbf{v}|\pi \mathbf{v})_{\mathcal{H}} = (P \pi \mathbf{v}|P \pi \mathbf{v})_{\mathcal{H}} = \|P \pi \mathbf{v}\|_{\mathcal{H}}^2 \leq \|\mathbf{v}\|_{\mathcal{H}}^2 = 1,$$

then $0 \leq \rho \leq 1$. Using (A.2) gives

$$|(P \mathcal{A} P \mathbf{v}|\mathbf{v})_{\mathcal{H}}| \leq \|h\|_{L^{\infty}(\Omega)} \|P \pi \mathbf{v}\|_{\mathcal{H}}^2.$$

It follows from (A.1) that $((Id - P)\mathcal{A}(Id - P)\mathbf{v}|\mathbf{v})_{\mathcal{H}} = 0$ and $\pi \mathcal{A} = \mathcal{A} \pi$, we have

$$(P \mathcal{A} P \mathbf{v}|\mathbf{v})_{\mathcal{H}} = (2\rho - 1)(\mathcal{A} \mathbf{v}|\mathbf{v})_{\mathcal{H}}$$

Since $|(\mathcal{A}\mathbf{v}|\mathbf{v})_{\mathcal{H}}| = |\sum_{i \in I} \lambda_i \frac{(\mathbf{v}|\phi_i)_{\mathcal{H}}^2}{\|\phi_i\|_{\mathcal{H}}^2}| \geq |\lambda_{-1}| \|\mathbf{v}\|_{\mathcal{H}}^2 = |\lambda_{-1}|$, we have

$$|\lambda_{-1}(2\rho - 1)| \leq \|h\|_{L^\infty(\Omega)} \|P\mathbf{v}\|_{\mathcal{H}}^2 = \|h\|_{L^\infty(\Omega)} \rho. \quad (\text{B.1})$$

This in turn implies that $\rho \geq \frac{|\lambda_{-1}|}{2|\lambda_{-1}| + \|h\|_{L^\infty(\Omega)}}$, hence there exists C independent of M such that

$$(\pi P \pi \phi | \phi)_{\mathcal{H}} \geq C \|\pi \phi\|_{\mathcal{H}} \quad \forall \phi \in \mathcal{H}(\mathcal{A}). \quad (\text{B.2})$$

Since $\pi_- \phi$ is the strong \mathcal{H} -limit of $\pi \phi$ as M goes to infinity and the constant C does not depend on M . Passing to the limit, we recover equation (B.2) with π replaced by π_- . The Lax-Milgram theorem applies and $\pi_- P \pi_-$ is a bijection from $\mathcal{R}(\pi_-)$ onto $\mathcal{R}(\pi_-)$ with a continuous inverse bounded by a constant in \mathcal{H} -norm.

We turn our interest to the bound in $1/2$ norm of B_- . Let $\phi \in \mathcal{D}(\mathcal{A}^{1/2})$, for any $M \in \mathbb{N}^*$ denote $\pi = \pi_{[-M, -1]}$, and let $\mathbf{v} = \pi B_- \phi$. We have $\pi \mathbf{v} = \mathbf{v}$ and $\mathbf{v} \in \mathcal{D}(\mathcal{A})$. Recalling (A.1) and $\pi \mathcal{A} = \mathcal{A} \pi$, we have

$$(P \mathcal{A} P \mathbf{v} | \mathbf{v})_{\mathcal{H}} = ((\mathcal{A} P + P \mathcal{A} - \mathcal{A}) \mathbf{v} | \mathbf{v})_{\mathcal{H}} = 2(P \mathbf{v}, \mathcal{A} \mathbf{v})_{\mathcal{H}} - (\mathcal{A} \mathbf{v}, \mathbf{v})_{\mathcal{H}}.$$

Hence, since $\pi \mathbf{v} = \mathbf{v}$ and π is a projection on negative eigenvalues of \mathcal{A} only, then $\|\mathbf{v}\|_{1/2}^2 = -(\mathcal{A} \mathbf{v} | \mathbf{v})_{\mathcal{H}}$ and

$$\|\mathbf{v}\|_{1/2}^2 = (P \mathcal{A} P \mathbf{v} | \mathbf{v})_{\mathcal{H}} - 2(P \mathbf{v} | \mathcal{A} \mathbf{v})_{\mathcal{H}} \leq \|h\|_{L^\infty(\Omega)} \|\mathbf{v}\|_{\mathcal{H}}^2 + 2\|\pi P \mathbf{v}\|_{1/2} \|\mathbf{v}\|_{1/2} \quad (\text{B.3})$$

Using the bound on the \mathcal{H} -norm of B_- , we have

$$\|\mathbf{v}\|_{\mathcal{H}} = \|\pi B_- \phi\|_{\mathcal{H}} \leq C \|\pi_- \phi\|_{\mathcal{H}} \leq C \|\pi_- \phi\|_{1/2}. \quad (\text{B.4})$$

We infer from (B.3) and (B.4) that $\|\mathbf{v}\|_{1/2} \leq C(\|\pi_- \phi\|_{1/2} + \|\pi P \mathbf{v}\|_{1/2})$. We let M go to infinity, then $\pi P \mathbf{v} = \pi P \pi B_- \phi$ goes to $\pi_- \phi$ and \mathbf{v} goes to $B_- \phi$, we obtain:

$$\|B_- \phi\|_{1/2} \leq C \|\pi_- \phi\|_{1/2},$$

which finishes the proof.

Appendix C. Proof of Theorem 5.

First case: $\mathcal{K}(\mathcal{A}) = \{0\}$, i.e. every case but the balanced case

In this case the condition $\pi_+ \psi = 0$ is then equivalent to $\psi = \pi_- \psi$.

After multiplication of (3.3) by $B_- \pi_-$, one obtains the following necessary condition for (3.3) to hold, which proves uniqueness:

$$\psi = B_- P \phi.$$

Denote $\theta = P B_- P \phi - P \phi$, the question of the existence of the solution is reduced to studying under which condition $\theta = 0$.

We have $P \theta = \theta$ and Theorem 4 states that $\pi_- \theta = 0$. This implies that $\theta \in \mathcal{R}(\pi_+)$. The operator \mathcal{A}^{-1} is symmetric definite positive on $\mathcal{R}(\pi_+)$, and induces the scalar product of the $-1/2$ -norm. Lemma 12 states that, if the constants are controlled we have $(\mathcal{A}^{-1} \theta | \theta)_{\mathcal{H}} = 0$, and it follows that $\theta = 0$. This proves the result when the constants are controlled.

Assume now that $\Gamma_D \cup \Gamma_R = \emptyset$ and $\int_{\Omega} h \neq 0$, Lemma 12 states that

$$\|\theta\|_{-1/2}^2 = (\mathcal{A}^{-1}\theta|\theta)_{\mathcal{H}} = \left(\int_{\Omega} h\right)^{-1} (\theta|\Phi)_{\mathcal{H}}^2. \quad (\text{C.1})$$

349 If $\int_{\Omega} h < 0$, the two terms have opposite signs, hence both are zero. Then $\theta = 0$ and
 350 this proves the result for the case $\Gamma_D \cup \Gamma_R = \emptyset$ and $\int_{\Omega} h < 0$.

Let us assume now that $\int_{\Omega} h > 0$. Since changing the sign of λ amounts to study the same problem where h is replaced by $-h$, we deduce from the case $\int_{\Omega} h < 0$ with $\phi = \Phi$ and the relation $P\Phi = \Phi$ that $PB_+\Phi = \Phi$. Since $\Phi \in \mathcal{D}(A^{1/2})$, it follows from Theorem 4 that $B_+\Phi \in \mathcal{D}(A^{1/2})$. Hence there exists a $s^* \in H$ such that $B_+\Phi = (1, s^*)$ and we have $\mathcal{A}B_+\Phi = (c^{-1}h - c^{-1}\operatorname{div} \sigma \nabla s^*, 0)$. This proves

$$P\mathcal{A}B_+\Phi = \mathcal{A}B_+\Phi,$$

and a simple calculation proves that

$$(\mathcal{A}B_+\Phi|\Phi)_{\mathcal{H}} = (\mathcal{A}B_+\Phi|B_+\Phi)_{\mathcal{H}} = \int_{\Omega} h \quad (\text{C.2})$$

We then compute

$$(\Phi - PB_-\Phi|\mathcal{A}B_+\Phi)_{\mathcal{H}} = (\Phi - B_-\Phi|P\mathcal{A}B_+\Phi)_{\mathcal{H}} = (\Phi - B_-\Phi|\mathcal{A}B_+\Phi)_{\mathcal{H}} \underbrace{=}_{(1)} (\Phi|\mathcal{A}B_+\Phi)_{\mathcal{H}} \neq 0,$$

351 where the equality (1) is obtained by remarking that $\mathcal{A}B_+\Phi \in \mathcal{R}(\pi_+)$ and $B_-\Phi \in$
 352 $\mathcal{R}(\pi_-)$ which are orthogonal spaces. We then obtain $\Phi - PB_-\Phi \neq 0$.

353 It follows from (C.1) that

$$\begin{aligned} \|\theta\|_{-1/2}^2 &= \left(\int_{\Omega} h\right)^{-1} (\theta|\Phi)_{\mathcal{H}}^2 = \left(\int_{\Omega} h\right)^{-1} (\theta|PB_+\Phi)_{\mathcal{H}}^2 \\ &= \left(\int_{\Omega} h\right)^{-1} (P\theta|B_+\Phi)_{\mathcal{H}}^2 = \left(\int_{\Omega} h\right)^{-1} (\theta|B_+\Phi)_{\mathcal{H}}^2. \end{aligned}$$

Using that θ and $B_+\Phi$ belong to $\mathcal{R}(\pi_+)$ on which all the eigenvalues of \mathcal{A}^{-1} are positive, the above equation implies

$$\|\theta\|_{-1/2}^2 = \left(\int_{\Omega} h\right)^{-1} (\theta|\mathcal{A}B_+\Phi)_{-1/2}^2.$$

We recall that $\|\mathcal{A}B_+\Phi\|_{-1/2}^2 = (\mathcal{A}B_+\Phi|B_+\Phi)_{\mathcal{H}} = \int_{\Omega} h$, and we obtain

$$\|\theta\|_{-1/2}^2 \|\mathcal{A}B_+\Phi\|_{-1/2}^2 = (\theta|\mathcal{A}B_+\Phi)_{-1/2}^2$$

which is an equality case in Cauchy-Schwarz inequality. This implies that θ and $\mathcal{A}B_+\Phi$ are colinear. Hence there exists some constant t such that

$$\theta = t\mathcal{A}B_+\Phi.$$

354 Performing the scalar product with Φ and using the fact that $(\mathcal{A}B_+\Phi|\Phi) \neq 0$
 355 which follows from (C.2), we conclude that $t = 0$ (hence $\theta = 0$) if and only if $(\theta|\Phi) = 0$,
 356 which reads $(\phi|\Phi - PB_-\Phi)_{\mathcal{H}} = 0$.

357 *Second case: $\mathcal{K}(\mathcal{A}) \neq \{0\}$, which is the balanced case.*

In the balanced case the kernel of \mathcal{A} is $\mathbb{R}\phi_0$, where we recall from section 2.1 that $P\phi_0 = \Phi$. The condition $\pi_+\psi = 0$ is equivalent to the existence of $\alpha \in \mathbb{R}$ such that $\psi = \pi_-\psi + \alpha\phi_0$. The condition $P\psi = P\phi$ is thus equivalent to

$$P\phi = P\pi_-\psi + \alpha\Phi. \quad (\text{C.3})$$

Necessary condition:

After multiplying Equation (C.3) by B_- one obtains:

$$\pi_-\psi = B_-P\phi - \alpha B_-\Phi.$$

Replacing the expression of $\pi_-\psi$ in (C.3) yields the following necessary condition:

$$PB_-P\phi + \alpha\Phi - \alpha PB_-\Phi = P\phi,$$

which reads

$$\alpha(\Phi - PB_-\Phi) = P\phi - PB_-P\phi.$$

It follows from lemma 13 that $(\Phi - PB_-\Phi|\Phi)_{\mathcal{H}} \neq 0$, then it is necessary that

$$\alpha = \frac{(\Phi - PB_-\Phi|\phi)_{\mathcal{H}}}{(\Phi - PB_-\Phi|\Phi)_{\mathcal{H}}}.$$

ψ is uniquely determined by

$$\psi = B_-P\phi + \frac{(\Phi - PB_-\Phi|\phi)_{\mathcal{H}}}{(\Phi - PB_-\Phi|\Phi)_{\mathcal{H}}}(\phi_0 - B_-\Phi) \quad (\text{C.4})$$

Conversely, if ψ is defined by (C.4), it is clear that $\pi_+\psi = 0$. Let $\theta = P\psi - P\phi$: it suffices to prove that $\theta = 0$ to ensure that ψ solves the problem.

$$(\theta|\phi_0)_{\mathcal{H}} = (\theta|\Phi)_{\mathcal{H}} = 0$$

by choice of α . A simple calculation shows that

$$\pi_-\theta = 0.$$

This proves that $\theta \in \mathcal{R}(\pi_+)$, where \mathcal{A}^{-1} is a symmetric positive definite operator. It follows from lemma 12 that $(\mathcal{A}^{-1}\theta|\theta)_{\mathcal{H}} = 0$ and hence $\theta = 0$. This finishes the proof.

REFERENCES

- [1] Gostoli C. and Gatta A. Mass transfer in a hollow fiber dialyzer. *Journal of Membrane Science*, 6:133–148, 1980.
- [2] Abram Dorfman and Zachary Renner. Conjugate problems in convective heat transfer. *Mathematical Problems in Engineering*, 2009, 2009.
- [3] MA Ebadian and HY Zhang. An exact solution of extended graetz problem with axial heat conduction. *International Journal of Heat and Mass Transfer*, 32(9):1709–1717, 1989.
- [4] Jérôme Fehrenbach, Frédéric de Gournay, Charles Pierre, and Franck Plouraboué. The generalized graetz problem in finite domains. *SIAM Journal on Applied Mathematics*, 72(1):99–123, 2012.
- [5] von L Graetz. Über die wärmeleitungsfähigkeit von flüssigkeiten. *Annalen der Physik*, 261(7):337–357, 1885.

- [6] Chii-Dong Ho, Ho-Ming Yeh, and Wen-Song Sheu. An analytical study of heat and mass transfer through a parallel-plate channel with recycle. *International journal of heat and mass transfer*, 41(17):2589–2599, 1998.
- [7] Kragh J., Rose J., T.R. Nielsen, and Svendsen S. New counter flow heat exchanger designed for ventilation systems in cold climates. *Energy and Buildings*, 39:1151–1158, 2007.
- [8] Ho-Eyoul Jeong and Jae-Tack Jeong. Extended graetz problem including streamwise conduction and viscous dissipation in microchannel. *International Journal of Heat and Mass Transfer*, 49(13):2151–2157, 2006.
- [9] J. Lahjomri, A. Oubarra, and A. Alemany. Heat transfer by laminar Hartmann flow in thermal entrance region with a step change in wall temperatures: the Graetz problem extended. *Int. J. Heat Mass Transfer*, 45(5):1127–1148, 2002.
- [10] ML Michelsen and John Villadsen. The graetz problem with axial heat conduction. *International Journal of Heat and Mass Transfer*, 17(11):1391–1402, 1974.
- [11] E. Papoutsakis, D. Ramkrishna, and H. C. Lim. The extended graetz problem with diriclet wall boundary conditions. *Appl. Sci. Res.*, 36:13–34, 1980.
- [12] E. Papoutsakis, D. Ramkrishna, and H. C. Lim. The extended graetz problem with prescribed wall flux. *AIChE J.*, 26:779–787, 1980.
- [13] E. Papoutsakis, D. Ramkrishna, and H-C. Lim. Conjugated graetz problems. pt.1: general formalism and a class of solid-fluid problems. *Chemical Engineering Science*, 36(8):1381–1391, 1981.
- [14] E. Papoutsakis, D. Ramkrishna, and H-C. Lim. Conjugated Graetz problems. Pt.2: Fluid-Fluid problem. *Chemical Engineering Science*, 36(8):1393–1399, 1981.
- [15] C. Pierre and F. Plouraboué. Numerical analysis of a new mixed-formulation for eigenvalue convection-diffusion problems. *SIAM Appl. Math.*, 70(3):658–676, 2009.
- [16] Dušan P. Sekulić R. K. Shah. *Fundamentals of heat exchanger design*. John Wiley and Sons, Ne Jersey, 2003.
- [17] Walter Rudin. *Functional analysis. International series in pure and applied mathematics*. McGraw-Hill, Inc., New York, 1991.
- [18] B Weigand, M Kanzamar, and H Beer. The extended graetz problem with piecewise constant wall heat flux for pipe and channel flows. *International journal of heat and mass transfer*, 44(20):3941–3952, 2001.